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Received March 28, 1978

(1415)

Invariant measures on the shift space

by

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Abstract. In this paper we investigate invariant measures on the space of sequences from a finite set S. Let p be an invariant measure on $X = \prod_{n=0}^{+\infty} S$ and let p_n be the joint distributions of p for $n = 1, 2, \ldots$ If p runs over all invariant measures on X, then the points p_n form a polygon K_n . We describe the set of all extremal points of K_n and we give a decomposition of Bernoulli measures by extremal points of K_n . Next, we study a class \mathcal{M}_0 of those measures which may be described by extremal points used in a decomposition of the Bernoulli measures. Further, we construct a complete system of invariants of the dynamical systems induced by the measures belonging to \mathcal{M}_0 .

1. Notations and definitions. Let $S = \{0, 1, ..., s-1\}, s \ge 2$, be a finite alphabet and let $X = \prod_{i=0}^{+\infty} S_i$. If $x = \{..., x_{-1}, x_0, x_1, ...\}$ is a point of X, then we define $T(x)_i = x_{i+1}$, $i = 0, \pm 1, \pm 2, ...$, that is, T shifts every sequence. Let # be a o-field of borelian subsets of X. A Borel probability measure p on B is called T-invariant (or shortly invariant) if $p(T^{-1}A) = p(A)$, for any $A \in \mathcal{B}$. For $n \ge 1$ we put $X_n = \prod_{i=1}^n S_i$. An element $B = (i_0 i_1 \dots i_{n-1})$ of X_n will be called a block. We shall identify Bwith the cylinder $\{x \in X; x_0 = i_0, x_1 = i_1, \ldots, x_{n-1} = i_{n-1}\}$. Let us denote by M(X) the set of all T-invariant measures on \mathcal{B} . For a given $p \in M(X)$ we define a measure p_n on X_n as $p_n(B) = p(B)$, $B \in X_n$, $n \ge 1$. The measure sure p_n may be considered as a point of the space R^{s^n} in the sense that the coordinates of p_n are indexed by the blocks $B \in X_n$, and the Bth coordinate of p_n is equal to $p_n(B)$. Fix $n \ge 1$ and denote by K_n the set of all vectors of the form $\langle p_n(B) \rangle_{B \in X_n}$, where p runs over all invariant measures on X. It is well known that the set K_n may be described by the following conditions:

$$\sum_{B\in X_n} p_n(B) = 1,$$

$$(b) \qquad \qquad \sum_{i=0}^{s-1} \, p_n(\mathit{C}i) = \sum_{i=0}^{s-1} \, p_n(i\mathit{C}), \quad \text{ for every } \mathit{C} \in \mathit{X}_{n-1},$$

$$(c) p_n(B) \geqslant 0, B \in X_n.$$

Further, if the measures p_n , n=1,2,..., are appointed by the invariant measure p, then the conditions of consistency are satisfied, i.e.

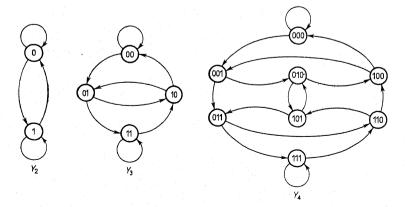
(d)
$$\sum_{i=0}^{s-1} p_{n+1}(Bi) = p_n(B), \quad B \in X_n, \ n = 1, 2, \dots$$

Condition (d) may be regarded as a definition of a mapping f_n from K_{n+1} onto K_n . We remark that the sets K_n , $n \ge 1$, are polygons in R^{s^n} and it is easy to check that $\dim K_n = s^{n-1}(s-1)$. We obtain a sequence of the polygons K_n and the functions f_n ,

$$K_1 \stackrel{f_1}{\leftarrow} K_2 \stackrel{f_2}{\leftarrow} K_3 \stackrel{f_3}{\leftarrow} \dots$$

In view of the above remarks the set M(X) may be identified with $\lim_{n \to \infty} K_n$. If $\overline{p} \in K_n$, $\overline{q} \in K_{n+1}$ and $\overline{p} = f_n(\overline{q})$, then we shall say that the vector \overline{q} is an extension of \overline{p} .

2. Extremal points of K_n . Now, we shall describe the set of all extremal points of K_n . In order to do this we use a graph Y_n , $n=2,3,\ldots$ If n=1, then K_1 may be identified with the simplex $T_s=\{(x_0,x_1,\ldots,x_{s-1}); \sum x_i=1, x_i \geq 0\}$, the extremal points of which are $(1,0,\ldots,0)$, $(0,1,0,\ldots,0),\ldots,(0,\ldots,0,1)$. The vertices of Y_n form the blocks $C \in X_{n-1}$ and two blocks $C_1=(i_0\ldots i_{n-2})$ and $C_2=(j_0\ldots j_{n-2})$ are joined by an oriented edge (write $(C_1,C_2)\in \mathscr{A}_n$) iff $(i_1\ldots i_{n-2})=(j_0\ldots j_{n-3})$. This means that the end of C_1 agrees with the beginning of C_2 . In the case of n=2 each two block-symbols are joined by edges. For example, if $S=\{1,0\}$ then Y_2,Y_3,Y_4 have the following form:



Observe that the edges of Y_n may be identified with the blocks of length n in the following sense: each edge (C_1, C_2) determines a block $B = (i_0, i_1, \dots, i_{n-2}, j_{n-2})$.

Let $\gamma = \{B_1, B_2, \ldots, B_l\}$, $B_i \in X_n$, $i = 1, 2, \ldots, l$, $1 \le l \le s^{n-1}$, be a closed path in Y_n not having any loop. Define a vector $\overline{p}_{\gamma} = \langle p_{\gamma}(B) \rangle_{B \in X_n}$ as follows:

$$p_{\gamma}(B) = \begin{cases} 1/l, & B \in \gamma, \\ 0, & B \notin \gamma. \end{cases}$$

It is easy to see that $\overline{p}_{\nu} \in K_n$. Now we can prove

THEOREM 1. A vector $\overline{p} \in K_n$ is an extremal point of K_n iff $\overline{p} = \overline{p}_{\gamma}$, where γ is a closed path in Y_n which does not contain any loop.

Proof. Sufficiency. Suppose that $\gamma = \{B_1, \ldots, B_l\}, \ 1 \leqslant l \leqslant s^{n-1},$ is a closed path without loops. Let $B_i = (b_0^i, b_1^i, \ldots, b_{n-1}^i), C_i = (b_0^i, \ldots, b_{n-2}^i),$ $i = 1, 2, \ldots, l$. The blocks C_1, C_2, \ldots, C_l are the vertices of γ and they are pairwise distinct since γ does not contain any loop. Further, the condition that $B_1, B_2, \ldots, B_l, B_1$ are the successive edges of γ implies $B_i = b_0^i C_{i+1}, \ i = 1, 2, \ldots, l-1, \ \text{and} \ B_l = b_0^l C_l.$ Assume $\overline{p}_{\gamma} = t \cdot \overline{p} + (1-t) \cdot \overline{q}$, where 0 < t < 1 and $\overline{p}, \overline{q} \in K_n$. Then p(B) > 0 implies $B \in \gamma$. Hence $p(C_l j) > 0$ implies $j = b_{n-1}^i$ for $i = 1, 2, \ldots, l$ and $p(j C_i) > 0$ mplies $j = b_0^{i-1}, i = 2, \ldots, l$, and $j = b_0^l$ for i = 1. In this way we obtain

$$p(B_1) = p(b_0^1 C_2) = \sum_{j=0}^{s-1} p(jC_2) = \sum_{j=0}^{s-1} p(C_2 j) = p(C_2 b_{n-1}^2) = p(B_2).$$

Similarly we can establish $p(B_2) = p(B_3) = \ldots = p(B_l)$. Therefore the condition $\sum_{i=1}^{l} p(B_i) = 1$ implies $p(B_i) = 1/l$, $i = 1, 2, \ldots, l$, i.e. $\overline{p} = \overline{P}_{\gamma}$. In the same manner we obtain $\overline{p}_{\gamma} = \overline{q}$, so that \overline{p}_{γ} is an extremal point of K_n .

Necessity. The polygon K_n is described by conditions (a), (b), (c). It is easy to remark that the order of the system of equations (a), (b) is equal to s^{n-1} . Take an extremal point $\overline{p} \in K_n$. It is well known that r $(r=s^n-s^{n-1})$ of the s^n coordinates of \overline{p} are equal to zero and the remaining s^{n-1} coordinates satisfy a regular subsystem of (a), (b). So p(B), $B \in X_n$, are rational numbers, say p(B) = r(B)/N, where r(B) are non-negative integers with $\sum_{B \in X_n} r(B) = N$. In order to find a closed path γ for which $\overline{p} = \overline{p}_{\gamma}$ we remark that condition (b) implies the following properties:

(1) for any $B \in X_n$ with p(B) > 0 there exists a $\overline{B} \in X_n$ such that $p(\overline{B}) > 0$ and $(B, \overline{B}) \in \mathscr{A}_{n+1}$.

73



Let $B_0 \in X_n$ be a block such that $p(B_0) = \min\{p(B): p(B) > 0\}$ Starting with B_0 and using (1), we may find a finite sequence of blocks B_0, B_1, \ldots of X_n such that $(B_i, B_{i+1}) \in \mathcal{A}_{n+1}$. Denote by C_i the block of length n-1 which forms the beginning of B_{i+1} and the end of B_i . It is clear that the sequence of blocks C_0, C_1, \ldots contains pairwise different blocks $C_{s+1}, C_{s+2}, \ldots C_m$, where s < m and $(C_m, C_{s+1}) \in \mathscr{A}_n$. Then the blocks B_{s+1}, \ldots, B_m form a closed path of Y_n without loops. Moreover, we have p(B) > 0 if $B \in \gamma$.

Let $m_0 \ge 1$ be the length of γ . Assume $\overline{p}_{\gamma} \ne \overline{p}$. Then $m_0 < N$ and therefore the vector $\bar{q} = \left(\bar{p} - \frac{m_0}{N} \; \bar{p}_{\gamma}\right) \cdot \frac{N}{N - m_0}$ is an element of K_n . Hence $\overline{p} = \frac{N - m_0}{N} \cdot \overline{q} + \frac{m_0}{N} \overline{p}_{\nu}$, which means that \overline{p} is not an extremal point of K_n . This leads us to a contradiction, so that the theorem is proved.

THEOREM 2. If \bar{p} is an extremal point of K_n , $n=1,2,\ldots$, then there exists exactly one $\overline{q} \in K_{n+1}$ such that $\overline{p} = f_n(\overline{q})$. Moreover, \overline{q} is an extremal point of K_{n+1} .

Proof. First we assume n=1. In this case the set of all extremal points of K_1 is identical with the set

$$\{\overline{p}_0,\,\overline{p}_1,\,\ldots,\,\overline{p}_{s-1}\}, \quad ext{where} \quad p_i(j) = egin{cases} 1, & i=j, \ 0, & i
eq j, \end{cases} \quad i,j \in S\,.$$

It is easy to check that the vectors

$$\overline{q}_0,\ldots,\overline{q}_{s-1}\in K_2, \quad \overline{q}_i(B)=egin{cases} 1, & B=(ii) \ 0, & B
eq (ii) \end{cases},$$

are the only vectors of K_2 such that $f_1(\bar{q}_i) = \bar{p}_i$.

Now, let $n \ge 2$ and take a path $\gamma = \{B_1, \ldots, B_l\}$ of Y_n not having loops. Then the vector $\overline{q} \in K_{n+1}$ is an extension of \overline{p}_n , iff the following conditions are satisfied [2]

(2)
$$\sum_{i=0}^{s-1} q(iCj) = p_{\gamma}(Cj), \quad j = 0, 1, ..., s-1, \\ \sum_{j=0}^{s-1} q(iCj) = p_{\gamma}(iC), \quad i = 0, 1, ..., s-1,$$

where C is any block of the length n-1. Therefore, in order to solve the systems of equations (2) it suffices to find matrices $Q(O) = \langle q(iOj) \rangle$, i, j = 0, 1, ..., s-1, satisfying (2) for every $C \in X_{n-1}$.

The following cases are possible:

- (i) $p_{\alpha}(iC) = 0$, $p_{\alpha}(Cj) = 0$ for every $i, j \in S$;
- (ii) there exists exactly one $i_0, j_0 \in S$ with $p_{\nu}(i_0C) = p_{\nu}(Cj_0) = 1/l$ and $p_{\nu}(iC) = 0 = p_{\nu}(Cj)$ for the remaining $i, j \in S$.

So the only solution of the system of equations (2) is a vector $\bar{q} \in K_{n+1}$ defined as follows: q(iCj) = 0, $i, j \in S$, if (i) holds and $q(i_0Cj_0) = 1/l$, $q(iCi) = 0, (i,j) \neq (i_0,j_0)$ if (ii) holds. It is easy to check that the vector \overline{q} is an extremal point of K_{n+1} appointed by the vertices $B_1, B_2, ..., B_l$. This completes the proof of the theorem.

Remark 1. Let $S = \{0, 1\}$ and let l_n denote the number of all extremal points of K_n . We can immediately check that $l_1 = 2$, $l_2 = 3$, $l_3 = 6$, $l_{*}=19$. At the same time the dimension of the sets $K_{1}, K_{2}, K_{3}, K_{4}$ is equal to 1, 2, 4, 8, respectively. This means that the polygons K_n cannot be simplexes for n=3,4.

§ 3. Decomposition of the Bernoulli measures. In this section we present a decomposition of Bernoulli measures by extremal points of $K_n, n = 1, 2, ...$

In order to do this we define a relation of φ -equivalence between the elements of X_n . The relation φ is defined as follows: two blocks $B, \overline{B} \in X_n \text{ are } q\text{-equivalent iff } \overline{B} = (i_t i_{t+1} \dots i_n i_1 \dots i_{t-1}) \text{ for some } 1 \leqslant t \leqslant n,$ where $B = (i_1 \dots i_n)$.

Denote by \mathcal{A}_n the set of all classes of the φ -equivalence. Observe that each element of \mathcal{B}_n contains at most n blocks of X_n . We shall show that if $\gamma \in \mathcal{B}_n$, then the edges of γ form a closed path of Y_n not having loops. It turns out that, for any Bernoulli measure p, the measures p_n may be described by the extremal points \overline{p}_{γ} , $\gamma \in \mathcal{B}_n$, n = 1, 2, ...

Suppose $\gamma = \{B_1, B_2, ..., B_l\}, l \leq n$. We may assume that if $B_1 = (i_1 i_2 \dots i_n)$ then $B_2 = (i_2 i_3 \dots i_n i_1)$, $B_3 = (i_3 \dots i_n i_1 i_2)$, and so on. Let $B_i = C_i c_n^i$, $C_i \in X_{n-1}$, $c_n^i \in S$. In order to show that γ forms a closed path of Y_n without loops it remains to prove that the blocks C_1, C_2, \ldots, C_l are pairwise distinct. First we observe that all blocks of γ have the same numbers of symbols. Now, if $C_1 = C_2$ then $c_n^1 \neq c_n^2$ because $B_1 \neq B_2$, and therefore the frequencies of the symbols in B_1 and B_2 are different. Thus $C_1 \neq C_2$, and similarly we obtain $C_i \neq C_j$ for $i \neq j$. This means that γ forms a closed path of Y_n not having any loop.

Let $n_0, n_1, \ldots, n_{s-1}$ be non-negative integers with $n_0 + \ldots + n_{s-1} = n$ and let $\mathcal{B}(n_0, n_1, \ldots, n_{s-1})$ be the set of all blocks of X_n containing the symbol 0 n_0 times, the symbol 1 n_1 times, and so on. Denote by b, the length of $\gamma, \gamma \in \mathcal{B}_n$. It is clear that

$$\sum_{\gamma \in \mathscr{B}(n_0, \dots, n_{s-1})} b_{\gamma} = \frac{n!}{n_0! \, n_1! \, \dots \, n_{s-1}!}.$$

Take a Bernoulli measure p on X given by a probability vector $\overline{q}=(q_0,q_1,\ldots,q_{s-1})$. It is not difficult to check that for $n=2,3,\ldots$ we have

$$p_n = \sum_{n_0 + \ldots + n_{s-1} = n} q_n^{n_0} \cdot q_1^{n_1} \ldots q_{s-1}^{n_{s-1}} \sum_{\gamma \in \mathscr{B}(n_0, \ldots, n_{s-1})} \overline{p}_{\gamma} \cdot b_{\gamma}.$$

Remark 2. For any block $B\in X_n$ there exists a unique class $\gamma\in \mathscr{Q}_n$ such that $B\in \gamma$. Accordingly the vectors \overline{p}_{γ} , $\gamma\in \mathscr{Q}_n$, are linearly independent and therefore they form a simplex L_n in K_n . In general, the dimension of L_n is smaller than $\dim K_n$.

In the sequel we denote by \mathcal{M}_0 the set of all invariant measures p on X for which $p_n \in L_n$ for n = 1, 2, ...

§ 4. Description of the class \mathcal{M}_0 . First we introduce the notation. If $B = (i_1 \dots i_l)$, $C = (j_1 \dots j_m)$ are two blocks, then we shall denote by BC the block $(i_1 \dots i_l j_1 \dots j_m)$. We start with the following

LEMMA 1. A measure p on X belongs to \mathcal{M}_0 iff for any two blocks B, C the condition

(3)
$$\sum_{i=0}^{s-1} p(BiC) = p(CB)$$

is satisfied.

Proof. Necessity. The condition $p \in \mathcal{M}_0$ implies that if $i \in S$ and B, C are two blocks then p(BiC) = p(CBi) and further

$$\sum_{i=0}^{s-1} p(BiC) = \sum_{i=0}^{s-1} p(CBi) = p(CB).$$

Sufficiency. Taking B or C as empty blocks, we find that p is an invariant measure on X. In order to prove that $p \in \mathcal{M}_0$ it suffices to show the equality p(iB) = p(Bi) for any symbol $i \in S$ and any block B. Using (3) we have

$$p(iB) = \sum_{j=0}^{s-1} p(Bji) = \sum_{j,k=0}^{s-1} p(ikBj) = \sum_{k=0}^{s-1} p(ikB) = p(Bi),$$

which completes the proof of the lemma.

DEFINITION 1. We say that an invariant measure p on X is symmetric if p(B) = p(C) for any blocks $B, C \in \mathcal{B}(n_0, n_1, \ldots, n_{s-1})$, where $n_0, n_1, \ldots, n_{s-1}$ are non-negative integers.

THEOREM 3. An invariant measure p on X belongs to \mathcal{M}_0 iff p is symmetric.

Proof. Sufficiency. If p is symmetric then p is constant on each class $\gamma \in \mathcal{B}_n$ for $n=1,2,\ldots$ But this means that for $n=1,2,\ldots,p_n \in L_n$, i.e. $p \in \mathcal{M}_0$.

Necessity. Let T be the shift on X and let k>1. Take blocks $B_0\in B_{l_0}$, $B_1\in X_{l_1},\ldots,B_k\in X_{l_k}$, where l_0,l_1,\ldots,l_k are positive integers and let $n_1>l_0$, $n_2>l_1+n_1,\ldots,n_k>n_{k-1}+l_{k-1}$. Then we have

$$\sum_{A_1, A_2, \dots, A_k} p(B_0 A_1 B_1 A_2 \dots A_k B_k)$$

$$= p(B_0 \cap T^{-n_1}(B_1) \cap T^{-n_2}(B_2) \cap \dots \cap T^{-n_k}(B_k)),$$

where $A_1 \in X_{n_1-l_0}$, $A_2 \in X_{n_2-n_1-l_1}$, $A_3 \in X_{n_3-n_2-l_2}$, and so on. Further, from the definition of \mathcal{M}_0 and (3) follows

$$\sum_{A_1,A_2,\ldots,A_k} p(B_0A_1B_1A_2\ldots A_kB_k) = p(B_0B_1\ldots B_k).$$

Thus for sufficiently large $n_k > n_{k-1} > \ldots > n_1$ we have

$$(4) p(B_0B_1...B_k) = p(B_0 \cap T^{-n_1}(B_1) \cap ... \cap T^{-n_k}(B_k)).$$

Now, take a partition ζ of X on ergodic components with respect to T. Let $M = X/\zeta$, $(p_m)_{m \in M}$ be conditional measures of ζ and let p_{ζ} be the quotient measure on M induced by p. Then we have

$$p(B_0 \cap T^{-n_1}(B_1) \cap \ldots \cap T^{-n_k}(B_k))$$

$$= \int_M p_m(B_0 \cap T^{-n_1}(B_1) \cap \ldots \cap T^{-n_k}(B_k)) p_{\xi}(dm).$$

By the above equality and by (4) we obtain

(5)
$$p(B_0B_1...B_k) = \int_M p_m(B_0 \cap T^{-n_1}(B_1) \cap ... \cap T^{-n_k}(B_k)) p_{\xi}(dm).$$

Now, fix $n_{k-1} > n_{k-2} > ... > n_1$. Applying the ergodic theorem to the dynamical systems $(X, \mathcal{B}, p_m, T), m \in M$, we have

(6)
$$\lim_{n_k \to \infty} \frac{1}{n_k} \sum_{l=0}^{n_k-1} p_m(B_0 \cap T^{-n_1}(B_1) \cap \dots \cap T^{-n_{k-1}}(B_{k-1}) \cap T^{-l}(B_k))$$

$$= p_m(B_k) p_m(B_0 \cap T^{-n_1}(B_1) \cap \dots \cap T^{-n_{k-1}}(B_{k-1})) \quad \text{for a.e. } m \in M.$$

Since $\frac{1}{n_k} \sum_{i=0}^{n_{k-1}} p_m(B_0 \cap \ldots \cap T^{-t}(B_k)) \leq 1$, $k \geq 1$, we can integrate both sides of (6) and we obtain

(7)
$$\lim_{n_k \to \infty} \frac{1}{n_k} \sum_{t=0}^{n_k-1} \int_{M} p_m(B_0 \cap \dots \cap T^{-n_{k-1}}(B_{k-1}) \cap T^{-t}(B_k)) = \int_{M} p_m(B_k) p_m(B_0 \cap \dots \cap T^{-n_{k-1}}(B_{k-1})) p_t(dm).$$

Further, (5) and (7) imply

(8)
$$p(B_k B_{k-1} \dots B_1 B_0)$$

= $\int_{\mathcal{D}} p_m(B_k) \cdot p_m(B_0 \cap T^{-n_1}(B_1) \cap \dots \cap T^{-n_{k-1}}(B_{k-1})) p_{\xi}(dm),$

for any sufficiently large integers $n_{k-1} > n_{k-2} > \ldots > n_1$. Repeating the above arguments for fixed $n_{k-2} > n_{k-3} > \ldots > n_1$ and for $n_{k-1} \to \infty$, we obtain

$$p(B_k B_{k-1} \dots B_1 B_0)$$

= $\int_{\mathcal{M}} p_m(B_k) p_m(B_{k-1}) p_m(B_0 \cap T^{-n_1}(B_1) \cap \dots \cap T^{-n_{k-2}}(B_{k-2})) p_{\xi}(dm).$

Proceeding in the same manner, we have

$$(9) p(B_k B_{k-1} \dots B_1 B_0) = \int_{\mathcal{M}} p_m(B_k) p_m(B_{k-1}) \dots (p_m) B_1 p_m(B_0) p_{\zeta}(dm),$$

for any blocks B_0, B_1, \ldots, B_k . Therefore, if B is a block having n_0 0's, n_1 1's, and so on, then (9) gives

(10)
$$p(B) = \int_{M} p_0(m) \cdot p_1(m) \dots p_{s-1}(m) p_{\zeta}(dm),$$

where $p_i(m) = p_m(\{i\}), i \in S$. Thus (10) implies the theorem.

Remark 3. If $p \in \mathcal{M}_0$ and p is an ergodic measure, then p_{ζ} is a δ -measure concentrated at a point $m_0 \in M$. Therefore p is a Bernoulli measure given by the probability vector $\langle p_0(m_0), \ldots, p_{s-1}(m_0) \rangle$.

EXAMPLE 1. Take a simplex T_s of the space R^{s-1} defined in the following way: $\overline{x}=(x_0,x_1,\ldots,x_{s-1})\in T_s$ iff $\sum\limits_{i=0}^{s-1}x_i=1$ and $x_i\geqslant 0,\,i=0,1,\ldots\ldots,s-1$. Assume that \overline{p} is a normalized, borelian measure on T_s . Then we can define a measure p on X as follows:

(11)
$$p(B) = \int_{T_s} x_0^{n_0} w_1^{n_1} \dots w_{s-1}^{n_{s-1}} \overline{p}(d\overline{x}),$$

where $B \in \mathcal{B}(n_0, n_1, \ldots, n_{s-1})$. It is easy to verify that p is an invariant measure on X and $p \in \mathcal{M}_0$.

Now, using the well-known theorem of de Finetti [6] we have the following

THEOREM 4. For every $p\in\mathcal{M}_0$ there exists a unique measure \overline{p} on T_s such that (11) holds.

§ 5. Isomorphism theorems. Consider a probability measure \overline{p} on T_s and let $p=\psi(\overline{p})$ be a measure on X defined by (11). The measure \overline{p} determines a dynamical system $Z(\overline{p})=\big(X,\mathscr{B},\psi(\overline{p}),T\big)$. In this section we

give a necessary and sufficient condition for two dynamical systems $Z(\overline{p}_1)$ and $Z(\overline{p}_2)$ to be isomorphic. First we describe a decomposition of $Z(\overline{p})$ on ergodic components.

Let \mathscr{M}_e be the set of all invariant ergodic measures on X. On \mathscr{M}_e we can define a topology [1] induced by neighbourhoods of the following form: $M(\varepsilon, \mu_0, C_1, \ldots, C_k) = \bigcap_{i=1}^k \{\mu \in \mathscr{M}_e; |\mu(C_i) - \mu_0(C_i)| < \varepsilon\}$, where $\mu_0 \in \mathscr{M}_0$, $\varepsilon > 0$ and C_1, C_2, \ldots, C_k are cylinders. Denote by \mathscr{B}_e the σ -field of borelian subsets of \mathscr{M}_e . Now, for a given probability measure r on \mathscr{B}_e , we can define a T-invariant measure $\tilde{\mu}$ on X by

(12)
$$\tilde{\mu}(C) = \int_{\mathcal{A}_{\sigma}} \mu(C) \nu(d\mu),$$

where G is a block. It is well known [4] that any invariant measure $\tilde{\mu}$ on X has the form (12). Moreover, the correspondence $v \to \tilde{\mu}$ is one-to-one. The measure v is called the decomposition of $\tilde{\mu}$ on ergodic components.

Now, we remark that the measure \overline{p} on T_s may be identified with a decomposition of $\psi(\overline{p})$ on ergodic components. In fact, the set T_s may be identified with the set \mathcal{M}_b of all Bernoulli measures on X. Moreover, the natural topology of T_s is identical with the restriction of the topology of \mathcal{M}_c to $\mathcal{M}_b(\mathcal{M}_b)$ is a closed subset of \mathcal{M}_c). Therefore, the measure \overline{p} may be regarded as a measure on \mathcal{M}_b and on \mathcal{M}_c , too. For this measure equation (12) reduces to (11).

Having a decomposition of $\psi(\overline{p})$ on ergodic components, we can construct a complete system of invariants of $Z(\overline{p})$. To do this we consider a function H on T_s defined by

$$H(\overline{x}) = -\sum_{i=0}^{s-1} x_i \log x_i.$$

Let $I_s = \langle 0, \log s \rangle$ and let ζ be a partition of T_s on the sets $C_a = \{\overline{x} \in T_s; H(\overline{x}) = a\}, \ a \in I_s$. The measure \overline{p} determines the quotient measure \hat{p} on $I_s = T_s/\zeta$ and conditional measures $\{p_a\}, \ a \in I_s$. Let $\{m_n(a)\}$ be the type of \overline{p}_a , that is, let $\{m_n\}$ be a sequence of measurable functions defined on I_s such that

$$\sum_{n=1}^{\infty} m_n(a) \leqslant 1, \quad m_{n+1}(a) \leqslant m_n(a), \quad m_n(a) \geqslant 0 \quad \text{ for } \quad n=1,2,\ldots,$$

and for almost all $a \in I_s$ with respect to \hat{p} . We obtain a pair $\theta(\overline{p}) = (\hat{p}, \{m_n(a)\}_{a \in I_s})$.

THEOREM 5. Given two probability measures \overline{p}_1 , \overline{p}_2 on T_s , the dynamical systems $Z(\overline{p}_1)$, $Z(\overline{p}_2)$ are isomorphic iff $\theta(\overline{p}_1) = \theta(\overline{p}_2)$.



The proof of the theorem can be obtained by using Ornstein's Isomorphism Theorem [3] and Roklin's decomposition theorem [5], which can be formulated in the following form:

THEOREM 6. Let $\tilde{\mu}_1, \tilde{\mu}_2$ be two invariant measures on X and let v_1, v_2 be the decomposition on ergodic components of $\tilde{\mu}_1$ and $\tilde{\mu}_2$, respectively. The dynamical systems $(X, \mathcal{B}, T, \tilde{\mu}_1)$ and $(X, \mathcal{B}, T, \tilde{\mu}_2)$ are isomorphic iff there exists an invertible measure-preserving transformation $S: (\mathcal{M}_e, v_1) \rightarrow (\mathcal{M}_e, v_2)$ such that for a.e. $(\text{mod } v_1) \ \mu \in \mathcal{M}_e$ the ergodic dynamical systems (X, \mathcal{B}, T, μ) and $(X, \mathcal{B}, T, S_{\mu})$ are isomorphic.

Remark 4. If s=2 then $m_1(\log 2)=1$, $m_n(\log 2)=0$ for $n\geqslant 2$ and for $a\in (0,\log 2)$ we must have $m_n(a)=0$ for $n=3,4,\ldots$ If s>2 then $m_1(\log s)=1$, $m_n(\log s)=0$, $n=2,3,\ldots$, and the remaining measures may have arbitrary types.

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Received November 7, 1977 Revised version May 24, 1978

(1364)

On certain subspaces of a nuclear power series spaces of finite type

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Dedicated to Professor

Arf on his 70th birthday

Abstract. Let $\alpha=(\alpha_n)$ denote a stable nuclear exponent sequence of finite type. It is shown that a subspace X of $\Lambda_1(a)$ with a basis is either isomorphic to a subspace of $\Lambda_\infty(a)$ or X has a complemented subspace which is isomorphic to a power series space of finite type. Also applications of this result to spaces of analytic functions are discussed.

Introduction. Throughout, we let $\alpha=(\alpha_n)$ denote a nuclear exponent sequence of finite type which is assumed to be stable [6] (i.e. (α_{2n}/α_n) is bounded). By subspace we mean a closed, infinite-dimensional subspace. Recently, Dubinsky [5] characterized Köthe spaces which are isomorphic to subspaces of a power series space $A_1(\alpha)$ of finite type. In particular, the power series space $A_{\infty}(\alpha)$ of infinite type is isomorphic to a subspace of $A_1(\alpha)$ ([3]). The main result of this note is the following.

THEOREM 1. A subspace X of $\Lambda_1(a)$ with a basis is either isomorphic to a subspace of $\Lambda_{\infty}(a)$ or X has a complemented subspace which is isomorphic to a power series space of finite type.

We note that a subspace of $\Lambda_{\infty}(a)$ cannot have a subspace isomorphic to a power series space of finite type ([20]). For the special case of $\alpha_n = n^{1/d}$, the corresponding power series space of infinite type is isomorphic to the space $O(C^d)$ of entire functions in d variables and the corresponding power series spaces of finite type is isomorphic to the space $O(\Delta^d)$ of functions analytic in the d-dimensional unit polycylinder ([15]). In the final section, we apply Theorem 1 to spaces of analytic functions and obtain the following result.

THEOREM 2. Let M be a Stein manifold of dimension d and assume that O(M) has a basis. Then O(M) is either isomorphic to $O(C^d)$ or O(M) has a complemented subspace isomorphic to a power series space of finite type.

We introduce some terminology in the following section, which leads to the proof of Theorem 1. For any undefined terminology we refer to [9], [12], and [6]. This research was supported by the Scientific and Technical Research Council of Turkey.